

Analysis of Computer Experiments with Functional Response

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Abstract

This paper is motivated by a computer experiment conducted for optimizing residual stresses in the machining of metals. Although kriging is widely used in the analysis of computer experiments, it cannot be easily applied to model the residual stresses because they are obtained as a profile. The high dimensionality caused by this functional response introduces severe computational challenges in kriging. It is well known that if the functional data are observed on a regular grid, the computations can be simplified using an application of Kronecker products. However, the case of irregular grid is quite complex. In this paper, we develop a Gibbs sampling-based expectation maximization algorithm, which converts the irregularly spaced data into a regular grid so that the Kronecker product-based approach can be employed for efficiently fitting a kriging model to the functional data.

KEY WORDS: EM Algorithm; Gaussian Process Model; Gibbs sampling; Kriging; Latin Hypercube Design; Optimization.

1. INTRODUCTION

Computer experiments (Santner et al., 2003; Fang et al., 2006) refer to those experiments that are performed using computers with the help of physical models and numerical methods, such as finite element analysis. Many computer experiment responses are collected in a functional form. That is, for each setting of the experiment, responses are collected over an interval of some index, such as space or time. The study of functional data is important because it can help us understand how the factors affect the shape of the resulting curve, which can have a bearing on the performance of the object or system under investigation.

This paper is motivated by a computer experiment in the machining of metals using a hard turning process. The main objective of the experiment is to study the residual stresses generated in the machined surface because they are known to influence the fatigue life and are also associated with distortion in machined parts. Nine machining process variables are considered in this simulation experiment and for each setting, the residual stress profiles over the depth at three different locations in the machined part are generated. Specifically, each residual stress profile measured in the tool feed rate direction contains 376 residual stress values output by the simulation code over a depth of 376 microns measured from the part surface (i.e. one residual stress value per micron). Figure 1 illustrates five sample residual stress profiles obtained under different machining process settings. We can see that the residual stress profiles vary with the settings. Therefore, it is important to predict the effect of the process variables on the residual stress profile so that the machining process can be optimized to enhance the fatigue life of machined components.

The literature on modeling computer experiments with functional responses is scant as most of the modeling techniques focus on single or multiple outputs (Conti et al. 2009, Conti and O’Hagan 2010). Kriging (Sacks et al. 1989) is the most popular technique for modeling these data due to its interpolating property, which is desirable in deterministic computer experiments. However, kriging is not used for analyzing functional data because of the computational problems caused by the high dimensionality of functional responses, especially when they are collected with an intensive sampling rate. Techniques such as

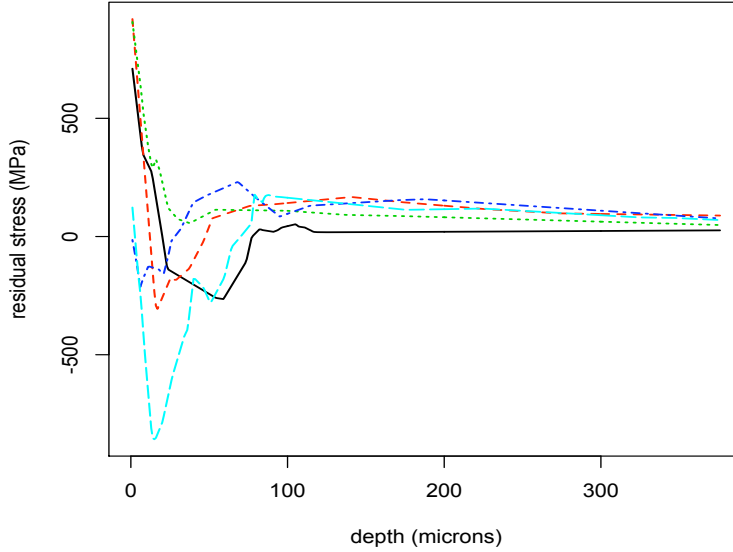


Figure 1: Five randomly sampled residual stress profiles vary with the settings.

wavelet decomposition (Bayarri et al. 2007), principal component analysis (PCA) (Ramsay and Silverman 2005, Higdon et al. 2007), functional linear models (Fang et al. 2006), and knot-based Gaussian process models (Banerjee et al. 2008) have been used instead. Although they provide simple and fast solutions, these models cannot interpolate the data. This creates a mismatch in the analysis methods for single and functional outputs. Thus, the main objective of this work is to improve the computational efficiency of kriging in the analysis of functional data so that the same method can be applied irrespective of the type of the response.

A naive extension of kriging to functional response is to include the functional argument as an additional input to the model (Kennedy and O’Hagan 2001, Liu and West 2009). For example, the machining experiments are conducted with nine process variables and a spatial variable that indicates the locations of the residual stress measurements. The kriging model can be fitted by incorporating depth as the 11th variable. Although this method is simple, it suffers from the computational difficulties. This is because the maximum likelihood estimation of the correlation parameters in the kriging model involves the inversion and

determinant calculation of a correlation matrix, whose dimension increases with the total number of observations $N = nm$, where n is the run size of the experiment, and m is the number of observations (from the functional space) in each run. These large matrix operations make the estimation of the correlation parameters computationally intensive and numerically unstable (Joseph and Kang 2011). Take the 90-run machining experiment as an example; the inversion and determinant calculation of a 33840×33840 matrix ($n = 90$, $m = 376$, and $N = 33840$) is required at each iteration of the optimization algorithm used in the maximum likelihood estimation, which will make the optimization extremely time-consuming.

One approach to overcome the computational issues associated with the naive kriging extension is to apply a Kronecker product formulation for constructing the correlation matrices (Williams et al. 2006, Rougier 2008, Liu et al. 2008, Bayarri et al. 2009). However, this approach can only be used when the functional responses are collected over a regular grid, which means that the outputs are observed at the same locations in the functional space for all runs.

Although not as common as the case of regular grid, observations on nonregular grid also occur sometimes in practice. For examples, in the study of transient rolling adhesion and deformation of leukocytes, the displacement profile of a cell is simulated using computational fluid dynamic models. These profiles are often truncated irregularly at various time points due to computational constraints at different kinetic parameter settings (Pappu and Bagchi 2008). In the thermo-mechanical study of the friction drilling processes, the thrust force profiles are generated from finite element modeling over the distance the tool travels. The profiles are irregularly collected because the travel distances are different at different values of tool feed rate (Miller and Shih 2007). In the motor engine simulations reported in Liu et al. (2008), the acceleration profiles are truncated at different time points leading to functional data on an irregular grid. Such truncated profiles are also commonly seen in degradation studies and fatigue testing simulations because the profile data cannot be collected when the product fails.

In this work, we propose a general and efficient method to overcome the computational

issue in analyzing functional response, especially with irregular grid. The paper is organized as follows. In Section 2, a brief review of the kriging model is given. The new modeling procedure for functional response in computer experiments is developed in Section 3. The proposed method is illustrated using the machining experiment in Section 4. Summary and concluding remarks are given in Section 5.

2. Kriging Preliminaries

Suppose that the computer experiment is conducted using p variables $\mathbf{x} = (x_1, \dots, x_p)'$ and the functional responses are collected over an index t . For each experimental setting $\mathbf{x}_i \in \mathcal{R}^p$, $i = 1, \dots, n$, the outputs $\mathbf{y}_i = (y_{i1}, \dots, y_{im_i})'$ are assumed to be a vector which is observed over t_{i1}, \dots, t_{im_i} with $\mathbf{t}_i = (t_{i1}, \dots, t_{im_i})'$. Note that the functional responses for each run can be collected differently in terms of location, therefore, \mathbf{t}_i 's and m_i 's are not necessarily the same for every i . Then, the kriging model is given by

$$y(\mathbf{x}, t) = \mathbf{v}(\mathbf{x}, t)' \boldsymbol{\mu} + Z(\mathbf{x}, t), \quad (1)$$

where $y(\mathbf{x}, t)$ is the response at point t and input setting \mathbf{x} , $\mathbf{v}(\mathbf{x}, t)' = (v_0(\mathbf{x}, t), v_1(\mathbf{x}, t), \dots, v_L(\mathbf{x}, t))$ is a set of known functions (usually, $v_0(\mathbf{x}, t) = 1$), and $\boldsymbol{\mu}$ is a vector of unknown parameters. We assume that $Z(\mathbf{x}, t)$ is a Gaussian process with mean 0 and the covariance function $\text{cov}\{y(\mathbf{x}_1, t_1), y(\mathbf{x}_2, t_2)\} = \sigma^2 r(\mathbf{x}_1 - \mathbf{x}_2, t_1 - t_2)$. Furthermore, it is common to assume a separable product correlation function: $r(\mathbf{x}_1 - \mathbf{x}_2, t_1 - t_2) = (\prod_{i=1}^p r_i(x_{i1} - x_{i2})) r_T(t_1 - t_2)$, where $r_i(x_{i1} - x_{i2})$ is the correlation function for the i th variable and $r_T(t_1 - t_2)$ is the correlation function for variable t . Unknown parameters associated with these correlation functions are denoted by $\boldsymbol{\xi}$. Note that the spatial-temporal model (Cressie 1993, Fang et al. 2006) is a special case of this model where \mathbf{x} represents space and t represents time.

Suppose the $N \times 1$ vector $\mathbf{y} = (\mathbf{y}'_1, \dots, \mathbf{y}'_n)'$ is a collection of all the functional outputs, where $N = \sum_i^n m_i$. The experimental design is denoted by the $N \times p$ matrix $\mathbf{X} = (\mathbf{1}'_{m_1} \otimes \mathbf{x}_1, \dots, \mathbf{1}'_{m_n} \otimes \mathbf{x}_n)' = (\mathbf{X}_1, \dots, \mathbf{X}_N)'$, where \otimes denotes the Kronecker product operator, $\mathbf{1}_{m_i}$ is a column of 1's having length m_i , and $\mathbf{T} = (\mathbf{t}'_1, \dots, \mathbf{t}'_n)' = (t_1^*, \dots, t_N^*)'$ is the corresponding

$N \times 1$ vector on the functional space. Based on model (1), the universal kriging predictor is given by

$$\hat{y}(\mathbf{x}, t) = \mathbf{v}(\mathbf{x}, t)' \hat{\boldsymbol{\mu}} + \mathbf{r}(\mathbf{x}, t)' \mathbf{R}_{\mathbf{x}, t}^{-1} (\mathbf{y} - \mathbf{V} \hat{\boldsymbol{\mu}}), \quad (2)$$

where $\mathbf{V} = (\mathbf{v}(\mathbf{X}_1, t_1^*), \dots, \mathbf{v}(\mathbf{X}_N, t_N^*))'$, $\hat{\boldsymbol{\mu}} = (\mathbf{V}' \mathbf{R}_{\mathbf{x}, t}^{-1} \mathbf{V})^{-1} \mathbf{V}' \mathbf{R}_{\mathbf{x}, t}^{-1} \mathbf{y}$, $\mathbf{r}(\mathbf{x}, t) = (r(\mathbf{x} - \mathbf{X}_1, t - t_1^*), \dots, r(\mathbf{x} - \mathbf{X}_N, t - t_N^*))'$, and $\mathbf{R}_{\mathbf{x}, t}$ is the $N \times N$ correlation matrix with elements $r(\mathbf{X}_i - \mathbf{X}_j, t_i^* - t_j^*)$. Note that, the kriging predictor in (2) interpolates the data because $\hat{y}(\mathbf{x}_i, t_{ij}) = \mathbf{v}(\mathbf{x}_i, t_{ij})' \hat{\boldsymbol{\mu}} + (y_{ij} - \mathbf{v}(\mathbf{x}_i, t_{ij})' \hat{\boldsymbol{\mu}}) = y_{ij}$. The correlation parameters are estimated by minimizing the negative log-likelihood (Santner et al. 2003, pp. 66)

$$\hat{\boldsymbol{\xi}} = \arg \min_{\boldsymbol{\xi}} \left[N \log \hat{\sigma}^2 + \log |\mathbf{R}_{\mathbf{x}, t}| \right], \quad (3)$$

where $\hat{\sigma}^2 = \frac{1}{N} (\mathbf{y} - \mathbf{V} \hat{\boldsymbol{\mu}})' \mathbf{R}_{\mathbf{x}, t}^{-1} (\mathbf{y} - \mathbf{V} \hat{\boldsymbol{\mu}})$. Other approaches, such as cross-validation and restricted maximum likelihood (Santner et al. 2003, Section 3.3) can also be applied.

As mentioned before, an optimization algorithm employed in (3) may require hundreds of evaluations of $\mathbf{R}_{\mathbf{x}, t}^{-1}$ and $|\mathbf{R}_{\mathbf{x}, t}|$, which makes the estimation of correlation parameters computationally prohibitive. To tackle the computational difficulty associated with functional response and successfully extend kriging, we propose a new procedure in the next section.

3. Analysis of Functional Response

3.1 The first stage

We propose to build the kriging model in two stages. In the first stage, we use the marginal profiles in \mathbf{x} and t to identify the functional form of the mean $\mathbf{v}(\mathbf{x}, t)$ and obtain the initial estimates of the correlation parameters. It can be described as follows. Consider the model

$$y(\mathbf{x}, t) = \mu_0 + \mathbf{k}'(t) \mathbf{u}_t + \mathbf{g}'(\mathbf{x}) \boldsymbol{\nu}_x + Z(\mathbf{x}, t), \quad (4)$$

where $\mathbf{k}(t) = (k_1(t), \dots, k_a(t))'$ and $\mathbf{g}(\mathbf{x}) = (g_1(\mathbf{x}), \dots, g_b(\mathbf{x}))'$. This is a special case of (1) because the mean does not entertain the possible interactions between \mathbf{x} and t . These interactions are now absorbed into $Z(\mathbf{x}, t)$.

First, we compute the average profile. Assume that the functional outputs are collected from m different locations of the functional space. If the functional data are collected over an irregular grid, m is defined to be the union of the abscissae of the individual profiles. We assume that there are n_j observations for location j , $j = 1, \dots, m$. For each j , define $\bar{e}_{\cdot j} = n_j^{-1} \sum_{i=1}^{n_j} (y_{ij} - \bar{y}_{i\cdot})$, where $\bar{y}_{i\cdot} = m_i^{-1} \sum_{l=1}^{m_i} y_{il}$. Note that we removed the effect of x -variables by subtracting $\bar{y}_{i\cdot}$ from y_{ij} . Based on $\bar{e}_{\cdot 1}, \dots, \bar{e}_{\cdot m}$, we fit the model

$$\bar{e}(t) = \mu_{t0} + \mathbf{k}'(t)\mathbf{u}_t + Z(t). \quad (5)$$

Next, we compute the average responses over the functional space, i.e., $\bar{y}_1, \dots, \bar{y}_n$, and fit the model

$$\bar{y}(\mathbf{x}) = \mu_{\mathbf{x}0} + \mathbf{g}'(\mathbf{x})\boldsymbol{\nu}_{\mathbf{x}} + Z(\mathbf{x}). \quad (6)$$

The unknown functions in $\mathbf{k}(t)$ and $\mathbf{g}(\mathbf{x})$ can be identified using a variable selection procedure, such as the Bayesian forward selection procedure implemented in blind kriging (Joseph et al. 2008). For the correlation functions, we consider the widely used power exponential correlation functions $r_i(x_{i1} - x_{i2}) = \exp\{-\alpha_i |x_{i1} - x_{i2}|^d\}$ for \mathbf{x} and $r_T(t_1 - t_2) = \exp\{-\beta |t_1 - t_2|^d\}$ for t , where $d = 1$ and $d = 2$ correspond to the exponential and Gaussian correlation functions, respectively. The estimates of the correlation parameters are denoted by $\boldsymbol{\alpha}^{(0)} = (\alpha_1^{(0)}, \dots, \alpha_p^{(0)})'$ and $\beta^{(0)}$, which are used as initial values in the second stage estimation.

The second stage, described in Sections 3.2 and 3.3, is implemented depending on how the functional outputs are collected.

3.2 The second stage: regular grid

Suppose the functional output is collected over a regular grid, i.e., $\mathbf{t}_1 = \dots = \mathbf{t}_n = \mathbf{t}$ and $m_1 = \dots = m_n = m$. Thus, in this situation, the functional outputs are observed at the same locations of the functional space for each experimental setting. It follows that $\mathbf{R}_{\mathbf{X}, \mathbf{t}} = \mathbf{R}_{\mathbf{X}} \otimes \mathbf{R}_{\mathbf{t}}$, where $\mathbf{R}_{\mathbf{X}}$ is a n -by- n correlation matrix corresponding to \mathbf{x} with elements $r(\mathbf{x}_i - \mathbf{x}_j)$ and $\mathbf{R}_{\mathbf{t}}$ is a m -by- m correlation matrix according to the \mathbf{t} components (Williams et al. 2006, Rougier 2008, Liu et al. 2008, Bayarri et al. 2009). We have $\mathbf{R}_{\mathbf{X}, \mathbf{t}}^{-1} = \mathbf{R}_{\mathbf{X}}^{-1} \otimes \mathbf{R}_{\mathbf{t}}^{-1}$. Consequently, the computational complexity involved in finding the inverse of $\mathbf{R}_{\mathbf{X}, \mathbf{t}}$ is dramatically reduced from $O(n^3 m^3)$ to $O(n^3 + m^3)$ (An and Owen 2001). The resulting

kriging predictor can be written as

$$\hat{y}(\mathbf{x}, t) = \mathbf{v}(\mathbf{x}, t)' \hat{\boldsymbol{\mu}} + \mathbf{r}(\mathbf{x}, t)' \left(\mathbf{R}_{\mathbf{X}}^{-1} \otimes \mathbf{R}_{\mathbf{t}}^{-1} \right) (\mathbf{y} - \mathbf{V} \hat{\boldsymbol{\mu}}), \quad (7)$$

where $\mathbf{v}(\mathbf{x}, t) = (1, \mathbf{k}'(t), \mathbf{g}'(\mathbf{x}))'$ and $\mathbf{V} = (\mathbf{v}(\mathbf{X}_1, t_1^*), \dots, \mathbf{v}(\mathbf{X}_N, t_N^*))'$. The maximum likelihood estimate (MLE) of the parameters are given by

$$\hat{\boldsymbol{\mu}} = \left[\mathbf{V}' \left(\hat{\mathbf{R}}_{\mathbf{X}}^{-1} \otimes \hat{\mathbf{R}}_{\mathbf{t}}^{-1} \right) \mathbf{V} \right]^{-1} \left[\mathbf{V}' \left(\hat{\mathbf{R}}_{\mathbf{X}}^{-1} \otimes \hat{\mathbf{R}}_{\mathbf{t}}^{-1} \right) \mathbf{y} \right], \quad (8)$$

$$\hat{\sigma}^2 = \frac{1}{N} (\mathbf{y} - \mathbf{V} \hat{\boldsymbol{\mu}})' \left(\hat{\mathbf{R}}_{\mathbf{X}}^{-1} \otimes \hat{\mathbf{R}}_{\mathbf{t}}^{-1} \right) (\mathbf{y} - \mathbf{V} \hat{\boldsymbol{\mu}}), \quad (9)$$

$$\hat{\boldsymbol{\xi}} = \arg \min_{\boldsymbol{\xi}} \left[N \log \hat{\sigma}^2 + m \log |\mathbf{R}_{\mathbf{X}}| + n \log |\mathbf{R}_{\mathbf{t}}| \right]. \quad (10)$$

The computational complexity can be further reduced using an idea in Finley et al. (2009), provided the grids are equally spaced. This is especially useful when the functional responses are collected with intensive sampling rate. Without loss of generality, assume that the equal spacing functional outputs are observed at locations $\mathbf{t} = (1, 2, \dots, m)$, where m can be extremely large. In this case, the evaluation of the kriging predictor can be simplified by using an exponential correlation function ($d = 1$). The use of exponential correlation function results in a non-differentiable predictor, however, lower degrees of smoothness can still be maintained for large m . The correlation matrix corresponding to \mathbf{t} can be written as

$$\mathbf{R}_{\mathbf{t}} = \begin{bmatrix} 1 & \rho & \rho^2 & \dots & \rho^m \\ \rho & 1 & \rho & \dots & \rho^{m-1} \\ \vdots & \vdots & \vdots & & \vdots \\ \rho^m & \rho^{m-1} & \rho^{m-2} & \dots & 1 \end{bmatrix}, \quad (11)$$

where $\rho = \exp(-\beta)$. It follows that $|\mathbf{R}_{\mathbf{t}}| = (1 - \rho^2)^m$ and

$$\mathbf{R}_{\mathbf{t}}^{-1} = \frac{1}{1 - \rho^2} \begin{bmatrix} 1 & -\rho & 0 & \dots & 0 \\ -\rho & 1 + \rho^2 & -\rho & \dots & 0 \\ \vdots & \vdots & \vdots & & \vdots \\ 0 & 0 & -\rho & 1 + \rho^2 & -\rho \\ 0 & 0 & 0 & -\rho & 1 \end{bmatrix}.$$

Because $\mathbf{R}_{\mathbf{t}}^{-1}$ is highly sparse and $|\mathbf{R}_{\mathbf{t}}|$ can be written in a simple closed form, the evaluation of kriging predictor (7) becomes much easier with a computational complexity of $O(n^3 + mn^2)$.

3.3 The second stage: irregular grid

Consider a more general situation where the functional responses are collected over an irregular grid. That is, for the i th experimental setting, functional response \mathbf{y}_i is observed at $\mathbf{t}_i = (t_{i1}, \dots, t_{im_i})'$, where \mathbf{t}_i 's and the lengths m_i 's are not necessarily the same for every i , $i = 1, \dots, n$. In this situation, the Kronecker product technique, which is the key to simplify the computation, cannot be directly utilized. We propose an intuitive idea to overcome the computational problem, which is to fill-in the “missing observations” so that the data appear to be on a regular grid and apply the procedure developed in the previous subsection. Even though the issue of missing data has been extensively studied in the literature, the idea of creating such missing observations to overcome the computational difficulty in kriging is new. Moreover, a direct application of the Expectation Maximization (EM) algorithm (Dempster et al. 1977) does not help in reducing the computational complexity. Here we introduce a carefully constructed iterative algorithm within the framework of EM algorithm to efficiently estimate the missing observations.

Define the missing data by $\mathbf{z} = (\mathbf{z}_1, \dots, \mathbf{z}_n)'$, where \mathbf{z}_i is a vector of missing observations in the i th run. Combining \mathbf{z} and the observed data $\mathbf{y} = (\mathbf{y}'_1, \dots, \mathbf{y}'_n)'$, we obtain the complete dataset $\mathbf{c} = (\mathbf{c}'_1, \dots, \mathbf{c}'_n)'$ with $\mathbf{c}_i = (\mathbf{y}'_i \cup \mathbf{z}'_i)'$, which can be viewed as the data collected on a regular grid $\mathbf{t} = (t_1, \dots, t_m)'$ with missing data \mathbf{z}_i on $\mathbf{t} \setminus \mathbf{t}_i$. The E-step in the EM algorithm is to obtain

$$Q(\boldsymbol{\theta} | \hat{\boldsymbol{\theta}}^{(k)}) = E(l_{\mathbf{c}}(\boldsymbol{\theta}) | \mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)}), \quad (12)$$

where

$$l_{\mathbf{c}}(\boldsymbol{\theta}) = \frac{N}{2} \log(2\pi) + \frac{N}{2} \log(\sigma^2) + \frac{1}{2} \log |\mathbf{R}_{\mathbf{x}, \mathbf{t}}| + \frac{1}{2\sigma^2} (\mathbf{c} - \mathbf{V}\boldsymbol{\mu})' \mathbf{R}_{\mathbf{x}, \mathbf{t}}^{-1} (\mathbf{c} - \mathbf{V}\boldsymbol{\mu}) \quad (13)$$

is the negative log-likelihood of the complete data and $\hat{\boldsymbol{\theta}}^{(k)} = (\hat{\boldsymbol{\mu}}^{(k)}, \hat{\sigma}^{2(k)}, \hat{\boldsymbol{\xi}}^{(k)})$ are the parameters estimated at the k th EM iteration. Clearly, the E-step in (12) requires the computation

of

$$E((\mathbf{c} - \mathbf{V}\boldsymbol{\mu})'\mathbf{R}_{\mathbf{X},\mathbf{t}}^{-1}(\mathbf{c} - \mathbf{V}\boldsymbol{\mu})|\mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)}) = (E(\mathbf{c}|\mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)}) - \mathbf{V}\boldsymbol{\mu})'\mathbf{R}_{\mathbf{X},\mathbf{t}}^{-1}(E(\mathbf{c}|\mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)}) - \mathbf{V}\boldsymbol{\mu}) + \text{tr}(\mathbf{R}_{\mathbf{X},\mathbf{t}}^{-1}\text{cov}(\mathbf{c}|\mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)})), \quad (14)$$

which cannot be done efficiently because the expressions of $E(\mathbf{c}|\mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)})$ and $\text{cov}(\mathbf{c}|\mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)})$ contain the inverse of the correlation matrix constructed using the observations on an irregular grid. This is the same computational hindrance we faced before. Therefore, such a direct application of EM cannot overcome the computational issue.

To simplify the computation, we modify the standard EM approach using an idea similar to Monte Carlo EM algorithm (Chan and Ledolter 1995). The main idea is to estimate the missing data run-by-run by conditioning on the complete data in the other runs. Let \mathbf{z}_i^j be the missing data estimated in the j th iteration and $\mathbf{c}^j = (\mathbf{c}_1^j, \dots, \mathbf{c}_n^j)$, where $\mathbf{c}_i^j = (\mathbf{y}_i \cup \mathbf{z}_i^j)$ are interspersed onto the regular grid \mathbf{t} . We can sample the missing data by using a Gibbs sampler according to the conditional distributions

$$\begin{aligned} \mathbf{c}_1^j &\sim f(\mathbf{c}_1|\mathbf{y}, \mathbf{z}_{(-1)}^{j-1}, \hat{\boldsymbol{\theta}}^{(j-1)}), \\ \mathbf{c}_2^j &\sim f(\mathbf{c}_2|\mathbf{y}, \mathbf{z}_1^j, \mathbf{z}_{(-1,-2)}^{j-1}, \hat{\boldsymbol{\theta}}^{(j-1)}), \\ &\vdots \\ \mathbf{c}_n^j &\sim f(\mathbf{c}_n|\mathbf{y}, \mathbf{z}_1^j, \dots, \mathbf{z}_{n-1}^j, \hat{\boldsymbol{\theta}}^{(j-1)}), \end{aligned} \quad (15)$$

with $\mathbf{z}_{(-i)}^j = (\mathbf{z}_1^j, \dots, \mathbf{z}_{i-1}^j, \mathbf{z}_{i+1}^j, \dots, \mathbf{z}_n^j)$. Because each step in (15) involves data from a regular grid, the computations are cheap. This is shown below.

The prior distribution (Currin et al. 1991) for \mathbf{c}_i is

$$f(\mathbf{c}_i|\hat{\boldsymbol{\theta}}^{(k)}) \sim N(\boldsymbol{\zeta}_{\mathbf{c}_i}, \boldsymbol{\Sigma}_{\mathbf{c}_i}), \quad (16)$$

where $\boldsymbol{\zeta}_{\mathbf{c}_i} = (\mathbf{v}(\mathbf{x}_i, t_1), \dots, \mathbf{v}(\mathbf{x}_i, t_m))'\boldsymbol{\mu}$ and $\boldsymbol{\Sigma}_{\mathbf{c}_i} = \sigma^2\mathbf{R}_{\mathbf{t}}$. It follows that

$$f(\mathbf{c}_i|\mathbf{y}_i, \hat{\boldsymbol{\theta}}^{(k)}) \sim N(\boldsymbol{\zeta}_i, \boldsymbol{\Sigma}_i), \quad (17)$$

$$f(\mathbf{c}_i|\mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}, \hat{\boldsymbol{\theta}}^{(k)}) \sim N(\boldsymbol{\zeta}_{(-i)}, \boldsymbol{\Sigma}_{(-i)}), \quad (18)$$

where

$$\begin{aligned}
\zeta_i &= \zeta_{\mathbf{c}_i} + (\mathbf{r}_{\mathbf{t}_i}(t_1), \dots, \mathbf{r}_{\mathbf{t}_i}(t_m))' \mathbf{R}_{\mathbf{t}_i}^{-1} [\mathbf{y}_i - (\mathbf{v}(\mathbf{x}_i, t_{i1}), \dots, \mathbf{v}(\mathbf{x}_i, t_{im_i}))' \boldsymbol{\mu}], \\
\boldsymbol{\Sigma}_i &= \boldsymbol{\Sigma}_{\mathbf{c}_i} - \sigma^2 (\mathbf{r}_{\mathbf{t}_i}(t_1), \dots, \mathbf{r}_{\mathbf{t}_i}(t_m))' \mathbf{R}_{\mathbf{t}_i}^{-1} (\mathbf{r}_{\mathbf{t}_i}(t_1), \dots, \mathbf{r}_{\mathbf{t}_i}(t_m)), \\
\zeta_{(-i)} &= \zeta_{\mathbf{c}_i} + (\mathbf{r}_{(-i)}(\mathbf{x}_i, t_1), \dots, \mathbf{r}_{(-i)}(\mathbf{x}_i, t_m))' \left(\mathbf{R}_{\mathbf{x}_{(-i)}}^{-1} \otimes \mathbf{R}_{\mathbf{t}}^{-1} \right) (\mathbf{c}_{(-i)} - \zeta_{\mathbf{c}_{(-i)}}), \\
\boldsymbol{\Sigma}_{(-i)} &= \boldsymbol{\Sigma}_{\mathbf{c}_i} - \sigma^2 (\mathbf{r}_{(-i)}(\mathbf{x}_i, t_1), \dots, \mathbf{r}_{(-i)}(\mathbf{x}_i, t_m))' \\
&\quad \left(\mathbf{R}_{\mathbf{x}_{(-i)}}^{-1} \otimes \mathbf{R}_{\mathbf{t}}^{-1} \right) (\mathbf{r}_{(-i)}(\mathbf{x}_i, t_1), \dots, \mathbf{r}_{(-i)}(\mathbf{x}_i, t_m)),
\end{aligned}$$

$\mathbf{r}_{\mathbf{t}_i}(t_j) = (r_T(t_j - t_{i1}), \dots, r_T(t_j - t_{im_i}))'$, $\zeta_{\mathbf{c}_{(-i)}} = E(\mathbf{c}_{(-i)})$, $\mathbf{R}_{\mathbf{x}_{(-i)}}$ is the $(n-1) \times (n-1)$ correlation matrix corresponding to \mathbf{x} except the i th setting, $\mathbf{y}_{(-i)}$ is the observed data except the i th profile, $\mathbf{z}_{(-i)}$ is the missing data except those from the i th run, and $\mathbf{c}_{(-i)} = (\mathbf{y}_{(-i)} \cup \mathbf{z}_{(-i)})$. Then, we obtain the following result.

Proposition 1: Based on (16), (17), and (18), it follows that

$$f(\mathbf{c}_i | \mathbf{y}, \mathbf{z}_{(-i)}, \hat{\boldsymbol{\theta}}^{(k)}) \sim N(\boldsymbol{\eta}_i, \boldsymbol{\Gamma}_i), \quad (19)$$

where

$$\boldsymbol{\eta}_i = (\boldsymbol{\Sigma}_{(-i)}^{-1} + \boldsymbol{\Sigma}_i^{-1} - \boldsymbol{\Sigma}_{\mathbf{c}_i}^{-1})^{-1} (\boldsymbol{\Sigma}_i^{-1} \zeta_i + \boldsymbol{\Sigma}_{(-i)}^{-1} \zeta_{(-i)} - \boldsymbol{\Sigma}_{\mathbf{c}_i}^{-1} \zeta_{\mathbf{c}_i}), \quad (20)$$

$$\boldsymbol{\Gamma}_i = (\boldsymbol{\Sigma}_{(-i)}^{-1} + \boldsymbol{\Sigma}_i^{-1} - \boldsymbol{\Sigma}_{\mathbf{c}_i}^{-1})^{-1}. \quad (21)$$

The proofs are given in the Appendix. Based on (19), the conditional mean of \mathbf{c}_i is a weighted average of three terms: the mean $\zeta_{\mathbf{c}_i}$ from prior distribution, the conditional mean based on the i th observed profile $\zeta_i = E(\mathbf{c}_i | \mathbf{y}_i, \hat{\boldsymbol{\theta}}^{(k)})$, and the conditional mean based on the complete data except the i th profile $\zeta_{(-i)} = E(\mathbf{c}_i | \mathbf{c}_{(-i)}, \hat{\boldsymbol{\theta}}^{(k)})$. The weights are determined by their corresponding variances, $\boldsymbol{\Sigma}_{\mathbf{c}_i}$, $\boldsymbol{\Sigma}_i$, and $\boldsymbol{\Sigma}_{(-i)}$, which are also components in the conditional variance. Evaluating the conditional distribution in (19) is computationally efficient because $\zeta_{\mathbf{c}_i}$ and ζ_i involve observations with size m and m_i which are easy to calculate. This is also true for the variances $\boldsymbol{\Sigma}_{\mathbf{c}_i}$ and $\boldsymbol{\Sigma}_i$. Furthermore, because $\mathbf{c}_{(-i)}$ is on a regular grid, $\zeta_{(-i)}$ and $\boldsymbol{\Sigma}_{(-i)}$ can be easily evaluated. By simple modifications, results in Proposition 1 can be extended to the conditional distributions in (15). For example, the posterior distribution $f(\mathbf{c}_i | \mathbf{y}, \mathbf{z}_1^j, \mathbf{z}_{(-1,-2)}^{j-1}, \hat{\boldsymbol{\theta}}^{(k)})$ can be obtained by replacing $\mathbf{z}_{(-i)}$ with the updated missing data $(\mathbf{z}_1^j, \mathbf{z}_{(-1,-2)}^{j-1})$ in (19).

After generating q samples using the Gibbs sampler, we can now approximate $Q(\boldsymbol{\theta}|\hat{\boldsymbol{\theta}}^{(k)})$ in (12) by

$$Q^{(q)}(\boldsymbol{\theta}|\hat{\boldsymbol{\theta}}^{(k)}) = \frac{1}{q} \sum_{j=1}^q l_{\mathbf{c}^j}(\boldsymbol{\theta}). \quad (22)$$

However, if the Gibbs sampler needs too many samples to attain convergence, the computational advantage of the foregoing procedure may diminish. Interestingly, we can further simplify the procedure and speed-up the convergence using the following idea. Instead of sampling from the conditional distributions, we can simply iterate on the conditional expectations as follows:

$$\begin{aligned} \mathbf{z}_1^j &= E(\mathbf{z}_1|\mathbf{y}, \mathbf{z}_{(-1)}^{j-1}, \hat{\boldsymbol{\theta}}^{(j-1)}), \\ \mathbf{z}_2^j &= E(\mathbf{z}_2|\mathbf{y}, \mathbf{z}_1^j, \mathbf{z}_{(-1,-2)}^{j-1}, \hat{\boldsymbol{\theta}}^{(j-1)}), \\ &\vdots \\ \mathbf{z}_n^j &= E(\mathbf{z}_n|\mathbf{y}, \mathbf{z}_1^j, \dots, \mathbf{z}_{n-1}^j, \hat{\boldsymbol{\theta}}^{(j-1)}), \end{aligned} \quad (23)$$

where $E(\mathbf{z}_i|\mathbf{y}, \mathbf{z}_{(-i)}, \hat{\boldsymbol{\theta}}^{(k)})$ represents the expectation based on the conditional distribution $f(\mathbf{c}_i|\mathbf{y}, \mathbf{z}_{(-i)}, \hat{\boldsymbol{\theta}}^{(k)})$. These conditional expectations can be easily computed using (20). It is worth noting that

$$\mathbf{R}_{\mathbf{X}} = \begin{bmatrix} \mathbf{R}_{\mathbf{X}_{(-i)}} & \boldsymbol{\varrho}_i \\ \boldsymbol{\varrho}_i' & 1 \end{bmatrix}$$

and

$$\mathbf{R}_{\mathbf{X}}^{-1} = \begin{bmatrix} \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} + \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} \boldsymbol{\varrho}_i \boldsymbol{\varrho}_i' \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} s_R^{-1} & -\mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} \boldsymbol{\varrho}_i s_R^{-1} \\ -\boldsymbol{\varrho}_i' \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} s_R^{-1} & s_R^{-1} \end{bmatrix} = \begin{bmatrix} \mathbf{A}_{(-i)} & \mathbf{a}_i \\ \mathbf{a}_i' & b_i \end{bmatrix}, \quad (24)$$

where $s_R = 1 - \mathbf{u}_i' \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} \boldsymbol{\varrho}_i$. Therefore, the evaluation of $\mathbf{R}_{\mathbf{X}_{(-i)}}^{-1}$ in (20) can be efficiently updated by

$$\mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} = \mathbf{A}_{(-i)} - \frac{\mathbf{a}_i \mathbf{a}_i'}{b_i}.$$

Because of this simplification, the complexity in estimating the missing observations is successfully reduced from $O((\sum_i^n m_i)^3)$ to $O(n^3 + \sum_i^n m_i^3 + m^3 + mn^2)$ per iteration and it can be further reduced to $O(n^3 + mn^2)$ by using the exponential correlation function (11) for the functional argument. The asymptotic convergence of \mathbf{z}_i^j 's follows from the proposition below.

Proposition 2: Assume that

$$\max_i \sum_{k=1, k \neq i}^n d_k \left\| \left[\mathbf{I}_{m \times m} + \boldsymbol{\Sigma}_{(-i)} (\boldsymbol{\Sigma}_i^{-1} - \boldsymbol{\Sigma}_{\mathbf{c}_i}^{-1}) \right]_{ik}^{-1} \right\|_2 < 1, \quad (25)$$

where $(d_1, \dots, d_{i-1}, d_{i+1}, \dots, d_n) = \mathbf{r}'_{(-i)}(\mathbf{x}_i) \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1}$, $\|\mathbf{D}\|_2 = \max_{\|\mathbf{s}\|_2=1} \|\mathbf{D}\mathbf{s}\|_2$, and $[\cdot]_{ik}$ indicates the block matrix with rows corresponding to the missing data at the i th run and columns corresponding to the missing data at the k th run. Then for all i , $\|\mathbf{z}_i^q - E(\mathbf{z}_i | \mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)})\|_2$ converges to $\mathbf{0}$ as $m \rightarrow \infty$ and $q \rightarrow \infty$.

The assumption in (25) can be interpreted intuitively. Recall that in Proposition 1 we have $\boldsymbol{\Sigma}_i = \boldsymbol{\Sigma}_{\mathbf{c}_i} - \sigma^2(\mathbf{r}_{\mathbf{t}_i}(t_1), \dots, \mathbf{r}_{\mathbf{t}_i}(t_m))' \mathbf{R}_{\mathbf{t}_i}^{-1}(\mathbf{r}_{\mathbf{t}_i}(t_1), \dots, \mathbf{r}_{\mathbf{t}_i}(t_m))$. When $m_i = 0$, the equality $\boldsymbol{\Sigma}_i = \boldsymbol{\Sigma}_{\mathbf{c}_i}$ holds; when $m_i \neq 0$, $|\boldsymbol{\Sigma}_{\mathbf{c}_i} - \boldsymbol{\Sigma}_i| > 0$ and it increases with m_i for fixed $\hat{\boldsymbol{\theta}}$. Similarly, it follows that $|\boldsymbol{\Sigma}_i^{-1} - \boldsymbol{\Sigma}_{\mathbf{c}_i}^{-1}| = |\boldsymbol{\Sigma}_{\mathbf{c}_i}^{-1}| \sigma^2(\mathbf{r}_{\mathbf{t}_i}(t_1), \dots, \mathbf{r}_{\mathbf{t}_i}(t_m))' \mathbf{R}_{\mathbf{t}_i}^{-1}(\mathbf{r}_{\mathbf{t}_i}(t_1), \dots, \mathbf{r}_{\mathbf{t}_i}(t_m)) |\boldsymbol{\Sigma}_i^{-1}|$ increases with m_i and consequently the left side of (25) decreases. Hence, this assumption suggests that the number of observations from the i th profile (m_i) should be reasonably large in order to accurately estimate the missing data.

The simplification made by (23) results in a new approximation to the E-step:

$$Q^{(q)}(\boldsymbol{\theta} | \hat{\boldsymbol{\theta}}^{(k)}) \approx l_{\mathbf{c}^q}(\boldsymbol{\theta}), \quad (26)$$

where $l_{\mathbf{c}^q}$ is the negative log-likelihood evaluated in (13) with \mathbf{c} replaced by $\mathbf{c}^q = (\mathbf{y}, \mathbf{z}_1^q, \dots, \mathbf{z}_n^q)$. This is equivalent to estimating the missing observations and then pretending they were known.

Now we can implement the M-step in the EM algorithm, which is to update the MLEs by

$$\hat{\boldsymbol{\theta}}^{(k+1)} = \operatorname{argmin}_{\boldsymbol{\theta}} Q^{(q)}(\boldsymbol{\theta} | \hat{\boldsymbol{\theta}}^{(k)}),$$

where $Q^{(q)}(\boldsymbol{\theta} | \hat{\boldsymbol{\theta}}^{(k)})$ is obtained from (26). Since the estimation herein is based on the complete data \mathbf{c}^q collected on regular grids, estimators according to (8), (9), and (10) can be computed efficiently.

There are some parameters that need to be pre-specified in the EM procedure. The initial estimates of $\hat{\boldsymbol{\theta}}^{(0)} = (\hat{\boldsymbol{\mu}}^{(0)}, \hat{\sigma}^{2(0)}, \hat{\boldsymbol{\alpha}}^{(0)}, \hat{\beta}^{(0)})$ can be obtained from the first stage of our model

building procedure. The initial estimates of the missing data can be obtained by adding the predictions from the two kriging models (5) and (6). Thus, the missing data of the i th profile, $\mathbf{z}_i^{(0)}$, at point t , is given by

$$\hat{y}(\mathbf{x}_i, t) = \hat{\mu}_{t0} + \hat{\mu}_{x0} + \mathbf{k}(t)' \hat{\mathbf{u}}_t + \mathbf{g}(\mathbf{x}_i)' \hat{\boldsymbol{\nu}}_x + \mathbf{r}_T(t)' \mathbf{R}_t^{-1} (\mathbf{e}_t - \mathbf{K} \hat{\mathbf{u}}_t - \hat{\mu}_{t0} \mathbf{1}_{m_i}) + \mathbf{r}(\mathbf{x}_i)' \mathbf{R}_X^{-1} (\bar{\mathbf{y}}_x - \mathbf{G} \hat{\boldsymbol{\nu}}_x - \hat{\mu}_{x0} \mathbf{1}_n),$$

where $\mathbf{K} = (\mathbf{k}(t_1), \dots, \mathbf{k}(t_m))'$, $\mathbf{G} = (\mathbf{g}(\mathbf{x}_1), \dots, \mathbf{g}(\mathbf{x}_n))'$, and \mathbf{R}_t^{-1} and \mathbf{R}_X^{-1} are evaluated at $\hat{\beta}^{(0)}$ and $\hat{\boldsymbol{\alpha}}^{(0)}$. The EM procedure is terminated when $\max_i |\hat{\theta}_i^{(k+1)} - \hat{\theta}_i^{(k)}| < \Delta$, where the tolerance Δ can be chosen based on the required accuracy.

Since the proposed approximation converges to the standard E-step when $q \rightarrow \infty$, the convergence of this procedure is guaranteed based on the results in the EM literature. More details can be found in Wu (1983), Chan and Ledolter (1995), and Fort and Moulines (2003).

4. RESIDUAL STRESS OPTIMIZATION

In this section, we revisit the machining experiment originally performed by Hung et al. (2009) and apply the proposed method to analyze the functional outputs. The objective of the experiment is to optimize a turning process for hardened bearing steel. Because of the computational challenge, Hung et al. (2009) optimized only the cutting forces, which is a single output problem, whereas the main objective of the experiment was to optimize the hard turning process with respect to residual stress, which is a functional output and is the focus of this paper.

Ten variables are considered in this experiment (Table 1). A 30-run orthogonal-maximin branching Latin hypercube design (BLHD) (Hung et al. 2009) is constructed based on the first nine variables and is given in Table 2. The cutting edge shape is labeled “1” and “2” to stand for chamfer and hone, respectively. The second and third factor, chamfer angle and length, are divided into 15 levels and the setting corresponding to the hone edge ($x_1 = 2$) is left blank because these two factors exist only when chamfer edge is considered. This is the main difference between BLHD and the traditional Latin hypercube design. The rest of the factors are set up at 30 levels. Since the cutting edge shape is a categorical variable, an isotropic correlation function (Hung et al. 2009) is used in the analysis. Exponential

Table 1: Factors and their ranges in the hard turning experiment

Factor	Ranges
x_1 : Cutting edge shape	chamfer or hone
x_2 : Chamfer Angle (degree)	17 \sim 20
x_3 : Chamfer Length (μm)	115 \sim 140
x_4 : Cutting edge radius (μm)	5 \sim 25
x_5 : Rake angle (degree)	$-15 \sim -5$
x_6 : Tool nose radius (mm)	0.4 \sim 1.6
x_7 : Cutting speed (m/min)	120 \sim 240
x_8 : Feed rate (mm/rev)	0.05 \sim 0.15
x_9 : Depth of cut (mm)	0.1 \sim 0.25
x_{10} : Location	1, 2, 3

correlation functions as defined in Section 3.1 are used for the rest of the factors. For each experimental setting, residual stresses are collected as a function of depth using a highly sophisticated finite element-based machining simulation software *AdvantEdge* (Third Wave Systems Inc., Minneapolis, MN). The simulations are computationally intensive which require 12 to 24 hours of running time for each experimental run.

The residual stress profiles are originally collected on a regular grid, but to illustrate the procedure, we artificially truncated them (randomly from 100 to 376 microns) so that the profiles are observed over an irregular grid. The assumption (25) holds in general with such a truncation and the left hand side of (25) decreases with respect to the percentage of truncation. Figure 2 illustrates five such profiles. The average of the 90 residual stress profiles is shown in Figure 3. Theoretically, it is expected that the residual stress will tend towards 0 as the depth increases. Therefore, a nonlinear model might fit this profile better. We fitted $S = \exp(-\lambda t)(\mu_0 + \mu_1 t + \mu_2 t^2)$ for the averaged residual stress profile S using least squares and obtained an estimate of λ as 0.005. Based on this result, we transformed the

Table 2: Orthogonal-maximin BLHD for the hard turning experiment

Run	x_1	x_2	x_3	x_4	x_5	x_6	x_7	x_8	x_9
1	1	1	6	15	23	7	9	18	10
2	1	2	11	25	3	25	14	25	19
3	1	3	3	4	20	18	18	5	26
4	1	4	14	9	6	6	27	7	17
5	1	5	8	16	8	21	2	2	1
6	1	6	1	17	10	5	25	19	25
7	1	7	12	29	26	15	5	14	12
8	1	8	5	26	16	30	22	15	6
9	1	9	15	7	13	26	7	11	27
10	1	10	10	1	29	20	23	6	5
11	1	11	2	20	21	27	10	20	29
12	1	12	7	8	11	14	4	29	21
13	1	13	13	22	9	1	24	27	9
14	1	14	4	10	2	24	28	13	13
15	1	15	9	28	25	13	17	3	28
16	2			19	5	9	1	8	20
17	2			14	28	17	6	21	24
18	2			6	17	4	16	12	4
19	2			11	1	12	15	4	8
20	2			27	22	8	30	24	16
21	2			21	14	23	19	10	22
22	2			23	18	22	12	28	3
23	2			3	27	3	3	26	14
24	2			13	15	19	29	16	30
25	2			24	12	2	11	1	18
26	2			18	24	28	8	17	2
27	2			12	30	11	26	9	11
28	2			2	4	16	13	30	15
29	2			30	7	10	20	23	7
30	2			5	19	29	21	22	23

original functional data by multiplying $\exp(\lambda t)$ and used them for the model fitting.

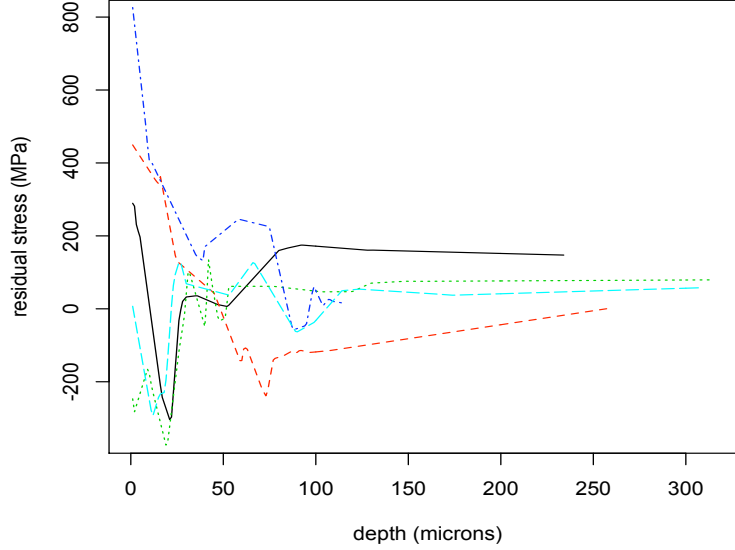


Figure 2: Truncated residual stress profiles.

By fitting $e(t) = \mu_{t0} + \mathbf{k}'(t)\mathbf{u}_t + Z(t)$ with $\mathbf{k}(t) = (t, t^2)'$, we obtained $\hat{\beta}^{(0)} = 1.24$. The blind kriging procedure (Joseph et al. 2008) applied on the 30 average values did not suggest any particular trend in the x -variables; therefore we fitted an ordinary kriging model $\bar{y}(\mathbf{x}) = \mu_{\mathbf{x}0} + Z(\mathbf{x})$ and obtained $\hat{\boldsymbol{\alpha}}^{(0)} = (0.10, 0.76, 0.45, 4.50, 3.72, 1.63, 0.17, 2.05, 4.50, 0.59)'$. Now we are ready to move on to the second stage of the model building procedure.

In the second stage, the proposed EM algorithm is applied. We applied 10 iterations ($q = 10$) using (23) inside each EM iteration. With $\Delta = 0.05$, the EM algorithm terminated after 21 iterations. Figure 4 shows some of the missing data estimated by the procedure. They are reasonably close to the true observations (the mean squared prediction error is 45.88), showing that the procedure works well. The maximum likelihood estimates of the correlation parameters are given by $\hat{\boldsymbol{\alpha}} = (1.34, 1.68, 1.47, 3.35, 2.80, 4.33, 1.82, 3.44, 4.09, 1.52)'$ and $\hat{\beta} = 1.12$. The computational saving obtained by the proposed EM procedure is quite substantial. The EM procedure took about 9.25 hours on a 3.4-GHz PC, whereas the naive kriging extension did not even converge after waiting for four days.

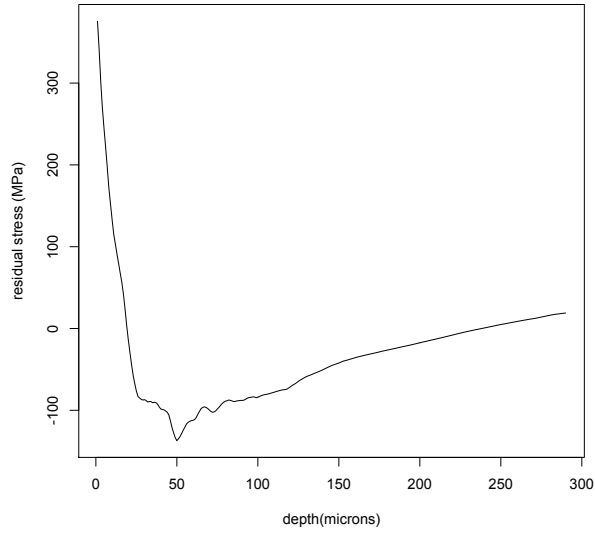


Figure 3: Average residual stress profile

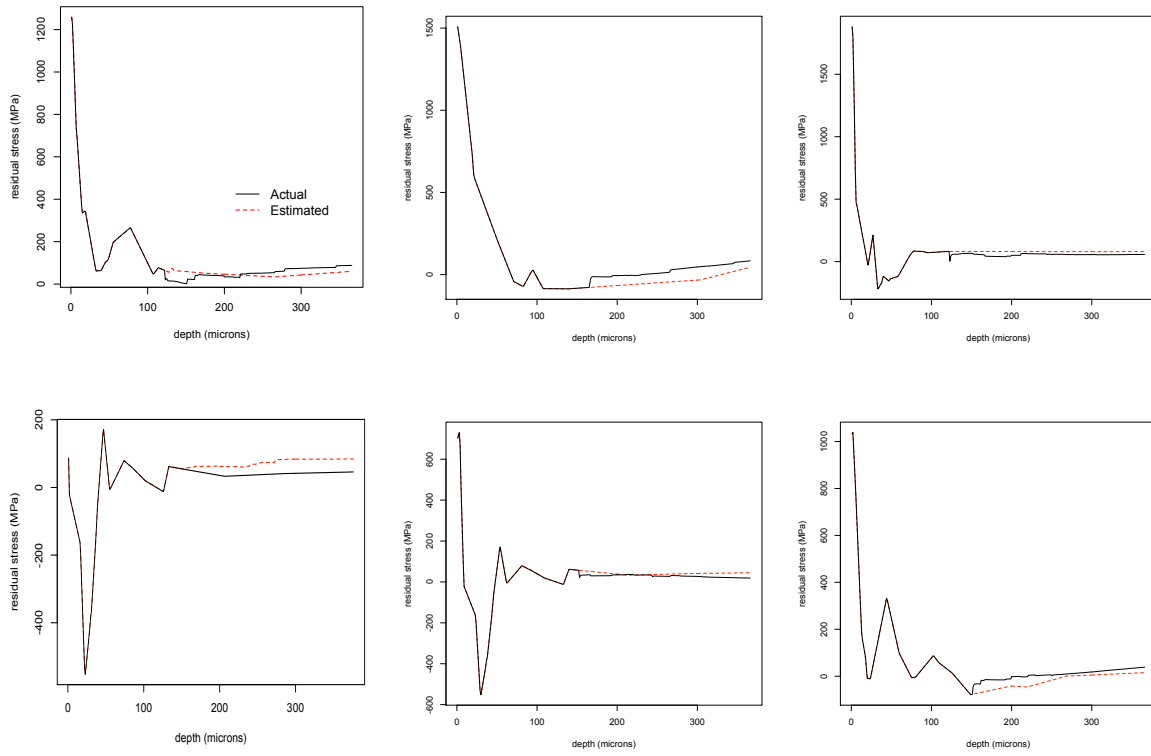


Figure 4: Missing data estimated by EM.

The $100(1 - \kappa)\%$ confidence interval for the prediction is given by (Santner et al. 2003, pp. 93),

$$\hat{y}(\mathbf{x}, t) \pm Z_{\kappa/2} \hat{\sigma} \{1 - \mathbf{r}(\mathbf{x}, t)' \mathbf{R}_{\mathbf{x}, t}^{-1} \mathbf{r}(\mathbf{x}, t) + \mathbf{h}' (\mathbf{V}' \mathbf{R}_{\mathbf{x}, t}^{-1} \mathbf{V})^{-1} \mathbf{h}\}^{1/2}, \quad (27)$$

where $Z_{\kappa/2}$ is the upper $\kappa/2$ critical point of the standard normal distribution and $\mathbf{h} = \mathbf{v}(\mathbf{x}, t) - \mathbf{V}' \mathbf{R}_{\mathbf{x}, t}^{-1} \mathbf{r}(\mathbf{x}, t)$. The 95% confidence intervals are illustrated using six randomly selected profiles (Figure 5). Each predicted curve was evaluated by leaving the original profile out from the training set. Note that we used a small nugget term to overcome the near-singularity in evaluating the large matrix inverse $\mathbf{R}_{\mathbf{x}, t}^{-1}$ (Gramacy and Lee 2011). The construction of the confidence intervals took about seven hours in a 2.33 GHz computer. This is not a big concern because we need to evaluate the matrix only once for constructing the confidence interval.

To illustrate the advantage of data augmentation, we compared the proposed approach with two approaches performed on the common grid (i.e., depths up to 100 microns). The two approaches include the widely used PCA (Higdon et al. 2007) and the naive kriging approach with functional argument as an additional input. In PCA, four basis functions were identified corresponding to the largest eigen values which account for more than 95% of the variation in the data. To compare the performance, we revisited the six randomly selected profiles and computed the mean squared cross validation error (MSCV) averaged over the entire depth. The MSCV for the proposed kriging method is 49.11, whereas the MSCV of the PCA method is 69.51 and the naive kriging method without data augmentation is 60.23. Although the PCA can be further improved by incorporating more basis functions, the improvement is marginal (e.g. with 10 basis functions, the MSCV is reduced only to 68.94.) The comparisons based on MSCV clearly show that the proposed kriging method provides better predictions than the other two methods. This is not surprising because without data augmentation, the information from 100 to 376 microns will be completely lost leading to severe extrapolation, which can be quite inaccurate.

The fitted model can be used for optimizing the residual stresses. The exact objective of optimization depends on the application area of the machined components. In most cases,

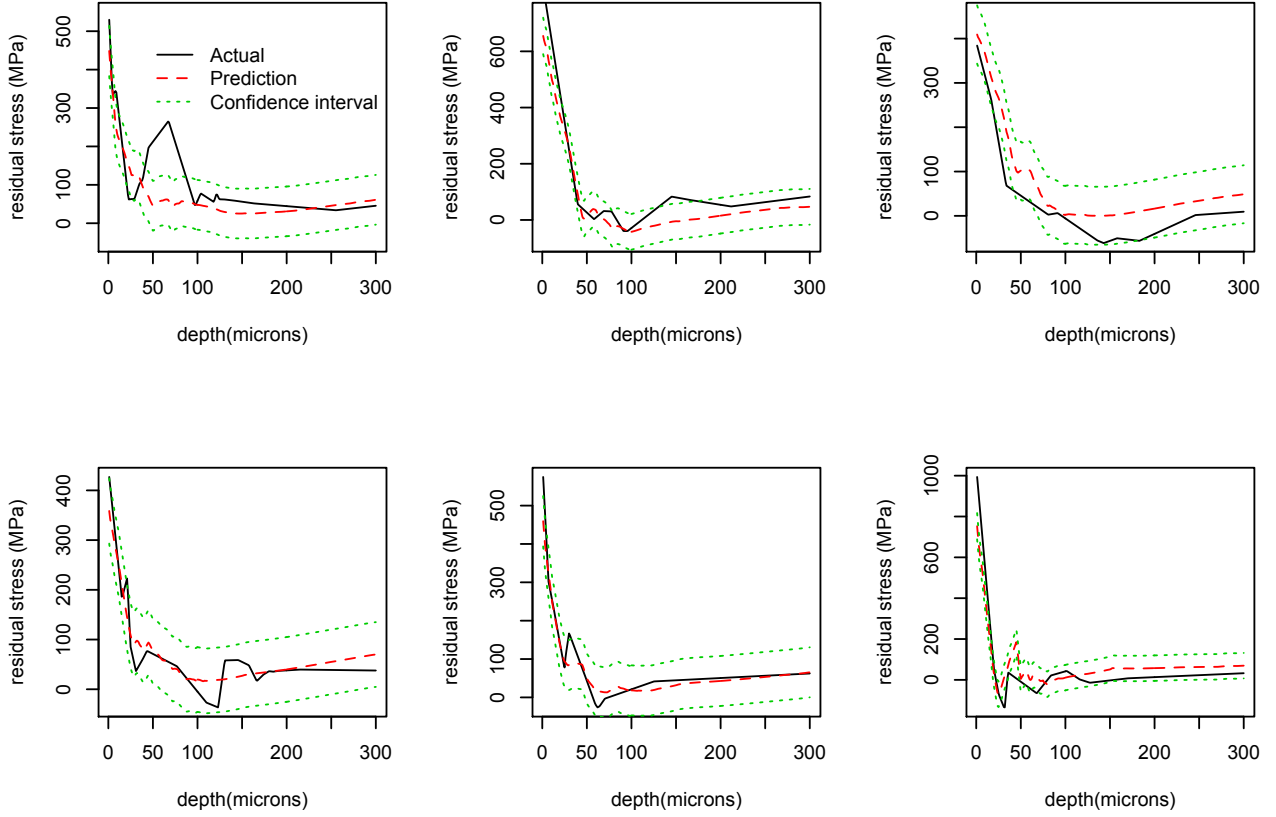


Figure 5: Leave-one-out prediction and confidence interval

the objective is to minimize the maximum tensile (i.e. maximum positive) residual stress. Thus, the objective is to find

$$\mathbf{x}^* = \min_{x_1, \dots, x_9} \left\{ \max_{t, x_{10}} \hat{y}(\mathbf{x}, t) \right\}. \quad (28)$$

The maximization over depth (t) and locations (x_{10}) provides the largest residual stress over the machined component. The optimal setting is given by $\mathbf{x}^* = (1.00, 11.38, 20.5, 11.49, 9.2, 6.14, 13.11, 3.99, 9.12)'$. Based on this setting, the predicted maximum residual stress is 58.84 MPa, which is about 53% smaller than the observed smallest maximum residual stress, 124.26 MPa, in the experiment. These results are reasonably close to that found using the original regular data, in which the optimal setting is $\mathbf{x}^* = (1.00, 11.61, 19.87, 11.44, 10.48, 7.66,$

12.63, 4.41, 9.35)' and the predicted maximum residual stress is 58.15 MPa. The predicted optimal residual stress profiles are illustrated in Figure 6 with the solid line indicating the prediction using the irregular data and the dash line indicating the prediction using the full regular data.

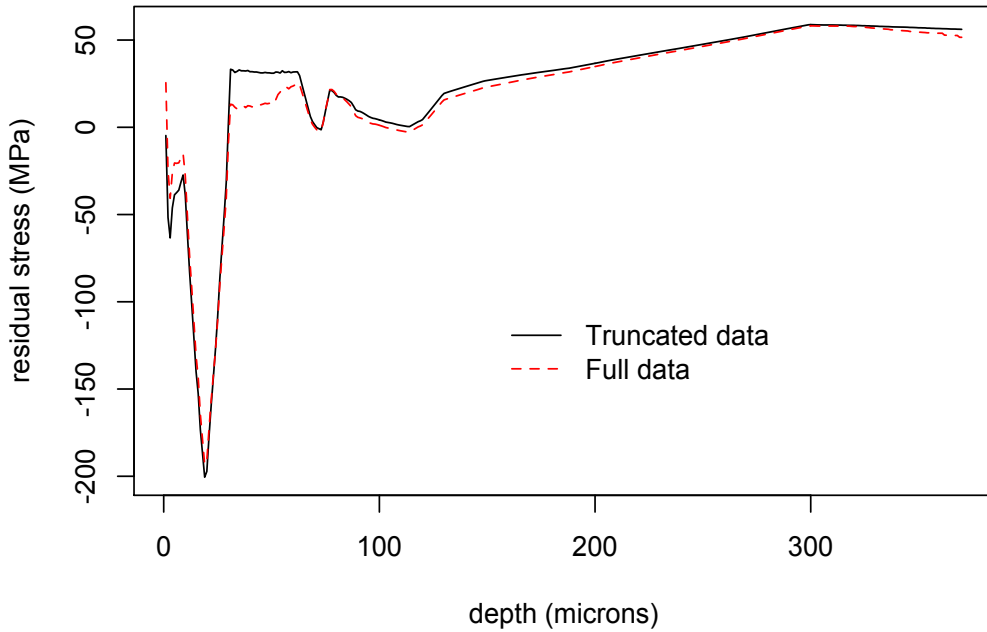


Figure 6: Optimal residual stress profiles with the full regular data and the truncated irregular data.

We also performed sensitivity analysis based on the main effect plots and two-factor interaction plots discussed in (Welch et al. 1992). The three main effects of feed, cutting edge shape, and the cutting speed, appear to be most important. Their effects on the residual stress profile are shown in Figure 7. For better clarity, the residual stress profiles were plotted with depths smaller than 100 microns (as the depth increases, the residual stresses converge to 0). We can see that the surface residual stress (i.e., depth=0) increases with feed. This trend is consistent with actual observations in the process. In Figure 7(b),

the hone edge tool leads to larger surface residual stress than with the chamfer edge tool. In Figure 7(c), the general effect of increasing the cutting speed is to cause the residual stress profile to become more tensile (or positive). This can be physically explained on the basis of the process physics and is attributed to increased thermal effects with increase in cutting speed.

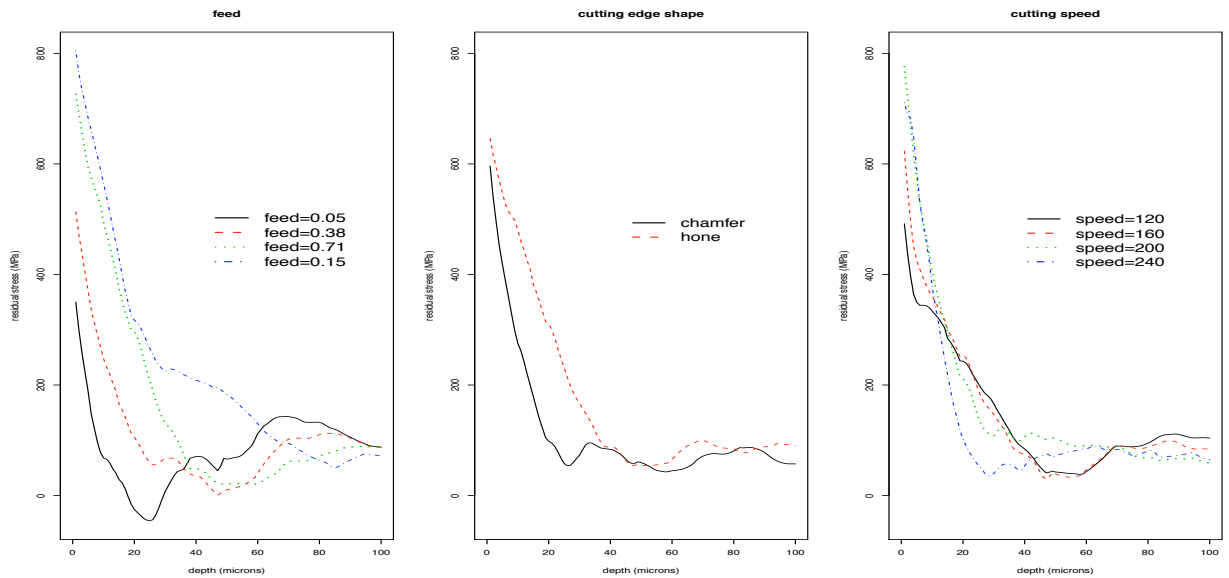


Figure 7: Significant main effects on the residual stress profile.

5. SUMMARY AND CONCLUDING REMARKS

Computer experiments with functional response are encountered in many engineering and scientific studies. Most available methods in the literature for analyzing computer experiments, however, focus on single outputs. Extending these methods to functional data analysis pose severe computational challenges. This article proposed a systematic methodology for modeling functional outputs using kriging especially for the data observed on an irregular grid. The computational challenge is successfully tackled by the proposed two-stage model building procedure which incorporates the Kronecker product technique and a version

of the EM algorithm for estimation. A direct application of the standard EM algorithm was not possible. Therefore, we introduced a Gibbs sampling-based computationally efficient algorithm for estimating the missing data run-by-run. The proposed method is illustrated by a hard turning machining simulation experiment in which the functional response, residual stress, is collected over depth. Based on the fitted model, the effects of each factor on the residual stress profile are studied and optimal settings are identified.

In this research, the uncertainty quantification is studied based on the plug-in approach, which underestimates the true variance. A possible extension is to develop a fully Bayesian approach by incorporating prior information from the mean coefficients and correlation parameters so that the uncertainty about the unknown parameters can be taken into account. However, a direct application of the fully Bayesian approach is computationally infeasible because each MCMC sample involves a large matrix inversion. Faster Bayesian computing techniques should be considered to tackle this problem (Joseph 2012).

A weakness of the current framework is the stationarity assumption employed in the kriging model. One approach to relax this assumption is to incorporate treed partitioning based on the idea of Chipman et al. (2002) and Gramacy and Lee (2008). Although additional computational efforts are required due to the recursive partitioning of trees, extension of this procedure is computationally tractable because the proposed algorithm can be implemented within each partition (leaf of the tree), which has less data.

Acknowledgments

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Appendix A: Proof of Proposition 1

The objective is to calculate conditional posterior $f(\mathbf{c}_i|\mathbf{y}, \mathbf{z}_{(-i)}) = f(\mathbf{c}_i|\mathbf{y}_i, \mathbf{y}_{(-i)}, \mathbf{z}_{(-i)})$.

We have

$$f(\mathbf{c}_i|\mathbf{y}_i, \mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}) \propto f(\mathbf{y}_i, \mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}|\mathbf{c}_i)f(\mathbf{c}_i),$$

and

$$f(\mathbf{y}_i, \mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}|\mathbf{c}_i) = f(\mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}|\mathbf{y}_i, \mathbf{c}_i)f(\mathbf{y}_i|\mathbf{c}_i).$$

Therefore,

$$\begin{aligned} f(\mathbf{c}_i|\mathbf{y}_i, \mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}) &\propto f(\mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}|\mathbf{y}_i, \mathbf{c}_i)f(\mathbf{y}_i|\mathbf{c}_i)f(\mathbf{c}_i) \\ &\propto f(\mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}|\mathbf{y}_i, \mathbf{c}_i)f(\mathbf{c}_i)f(\mathbf{y}_i|\mathbf{c}_i)f(\mathbf{c}_i)/f(\mathbf{c}_i). \end{aligned} \quad (29)$$

Because $\mathbf{y}_i \subset \mathbf{c}_i$, we have $f(\mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}|\mathbf{y}_i, \mathbf{c}_i) = f(\mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}|\mathbf{c}_i)$ and equation (29) can be written as

$$\begin{aligned} f(\mathbf{c}_i|\mathbf{y}_i, \mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}) &\propto f(\mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}|\mathbf{c}_i)f(\mathbf{c}_i)f(\mathbf{y}_i|\mathbf{c}_i)f(\mathbf{c}_i)/f(\mathbf{c}_i) \\ &\propto f(\mathbf{c}_i|\mathbf{y}_{(-i)}, \mathbf{z}_{(-i)})f(\mathbf{c}_i|\mathbf{y}_i)/f(\mathbf{c}_i). \end{aligned} \quad (30)$$

Since the three prior distributions on the right hand side of (30) are all normally distributed as given in (16), (17), and (18), we have

$$\begin{aligned} f(\mathbf{c}_i|\mathbf{y}_i, \mathbf{y}_{(-i)}, \mathbf{z}_{(-i)}) &\propto -\frac{1}{2} \exp \left\{ (\mathbf{c}_i - \boldsymbol{\zeta}_{(-i)})' \boldsymbol{\Sigma}_{(-i)}^{-1} (\mathbf{c}_i - \boldsymbol{\zeta}_{(-i)}) + (\mathbf{c}_i - \boldsymbol{\zeta}_i)' \boldsymbol{\Sigma}_i^{-1} (\mathbf{c}_i - \boldsymbol{\zeta}_i) \right. \\ &\quad \left. - (\mathbf{c}_i - \boldsymbol{\zeta}_{\mathbf{c}_i})' \boldsymbol{\Sigma}_{\mathbf{c}_i}^{-1} (\mathbf{c}_i - \boldsymbol{\zeta}_{\mathbf{c}_i}) \right\} \\ &\propto -\frac{1}{2} \exp \left\{ (\mathbf{c}_i - \boldsymbol{\eta}_i)' \boldsymbol{\Gamma}_i^{-1} (\mathbf{c}_i - \boldsymbol{\eta}_i) \right\}, \end{aligned}$$

where $\boldsymbol{\eta}_i$ and $\boldsymbol{\Gamma}_i$ are given in (20) and (21). Therefore, the conditional distribution of $f(\mathbf{c}_i|\mathbf{y}_i, \mathbf{y}_{(-i)}, \mathbf{z}_{(-i)})$ follows. \diamond

Appendix B: Proof of Proposition 2

First, it holds that

$$\begin{aligned}
& (\mathbf{r}_{(-i)}(\mathbf{x}_i, t_1), \dots, \mathbf{r}_{(-i)}(\mathbf{x}_i, t_m))' \left(\mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} \otimes \mathbf{R}_{\mathbf{t}}^{-1} \right) \\
&= \left(\mathbf{r}'_{(-i)}(\mathbf{x}_i) \otimes \mathbf{R}_{\mathbf{t}} \right) \left(\mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} \otimes \mathbf{R}_{\mathbf{t}}^{-1} \right) \\
&= \left(\mathbf{r}'_{(-i)}(\mathbf{x}_i) \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} \right) \otimes (\mathbf{R}_{\mathbf{t}} \mathbf{R}_{\mathbf{t}}^{-1}) \\
&= \left(\mathbf{r}'_{(-i)}(\mathbf{x}_i) \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1} \right) \otimes \mathbf{I}_{m \times m}.
\end{aligned}$$

By taking a factor $\Sigma_{(-i)}^{-1}$ out of (20) and assuming that $(d_1, \dots, d_{i-1}, d_{i+1}, \dots, d_n) = \mathbf{r}'_{(-i)}(\mathbf{x}_i) \mathbf{R}_{\mathbf{X}_{(-i)}}^{-1}$, we have

$$\begin{aligned}
\mathbf{z}_i^j &= \sum_{k=1}^{i-1} d_k \left[(\mathbf{I}_{m \times m} + \Sigma_{(-i)} \Sigma_i^{-1} - \Sigma_{(-i)} \Sigma_{\mathbf{c}_i}^{-1})^{-1} \right]_{ik} \mathbf{z}_k^j \\
&\quad + \sum_{k=i+1}^n d_k \left[(\mathbf{I}_{m \times m} + \Sigma_{(-i)} \Sigma_i^{-1} - \Sigma_{(-i)} \Sigma_{\mathbf{c}_i}^{-1})^{-1} \right]_{ik} \mathbf{z}_k^{j-1} + C,
\end{aligned} \tag{31}$$

where C is independent of \mathbf{z}_i 's. For notational simplicity, we will denote $[\cdot]_{ik}$ by $[\cdot]$ in the remaining part of the proof.

Consider a set of \mathbf{z}_i^j that satisfy (23) simultaneously, i.e., $\mathbf{z}_i^j \rightarrow \mathbf{z}_i$, we have

$$\begin{aligned}
\mathbf{z}_i &= \sum_{k=1}^{i-1} d_k \left[(\mathbf{I}_{m \times m} + \Sigma_{(-i)} \Sigma_i^{-1} - \Sigma_{(-i)} \Sigma_{\mathbf{c}_i}^{-1})^{-1} \right] \mathbf{z}_k \\
&\quad + \sum_{k=i+1}^n d_k \left[(\mathbf{I}_{m \times m} + \Sigma_{(-i)} \Sigma_i^{-1} - \Sigma_{(-i)} \Sigma_{\mathbf{c}_i}^{-1})^{-1} \right] \mathbf{z}_k + C.
\end{aligned} \tag{32}$$

Subtracting (32) from (31) and assuming that $\boldsymbol{\varepsilon}_i^j = \mathbf{z}_i^j - \mathbf{z}_i$, we obtain

$$\begin{aligned}
\boldsymbol{\varepsilon}_i^j &= \sum_{k=1}^{i-1} d_k \left[(\mathbf{I}_{m \times m} + \Sigma_{(-i)} \Sigma_i^{-1} - \Sigma_{(-i)} \Sigma_{\mathbf{c}_i}^{-1})^{-1} \right] \boldsymbol{\varepsilon}_k^j \\
&\quad + \sum_{k=i+1}^n d_k \left[(\mathbf{I}_{m \times m} + \Sigma_{(-i)} \Sigma_i^{-1} - \Sigma_{(-i)} \Sigma_{\mathbf{c}_i}^{-1})^{-1} \right] \boldsymbol{\varepsilon}_k^{j-1}.
\end{aligned}$$

Define $\delta = \max_i \sum_{k=1, k \neq i}^n d_k \left\| \left[\mathbf{I}_{m \times m} + \Sigma_{(-i)} \Sigma_i^{-1} - \Sigma_{(-i)} \Sigma_{\mathbf{c}_i}^{-1} \right]^{-1} \right\|_2$. Now we show by induction on i that the following result holds:

$$\max_i \|\boldsymbol{\varepsilon}_i^j\| \leq \delta \max_i \|\boldsymbol{\varepsilon}_i^{j-1}\|, \quad j = 1, 2, \dots \tag{33}$$

For $i = 1$,

$$\boldsymbol{\varepsilon}_1^j = \sum_{k=i+1}^n d_k \left[\mathbf{I}_{m \times m} + \Sigma_{(-i)} \Sigma_i^{-1} - \Sigma_{(-i)} \Sigma_{\mathbf{c}_i}^{-1} \right]^{-1} \boldsymbol{\varepsilon}_k^{j-1},$$

and thus, (33) clearly holds. Assume that (33) holds for $i = 1, \dots, l-1$ and based on assumption (25), we have for $i = l$

$$\begin{aligned}
\|\boldsymbol{\epsilon}_l^j\| &\leq \delta \max_l \|\boldsymbol{\epsilon}_l^{j-1}\| \left\| \sum_{k=1}^{l-1} d_k \left\| \left[(\mathbf{I}_{m \times m} + \boldsymbol{\Sigma}_{(-l)} \boldsymbol{\Sigma}_l^{-1} - \boldsymbol{\Sigma}_{(-l)} \boldsymbol{\Sigma}_{\mathbf{c}_l}^{-1})^{-1} \right] \right\|_2 \right. \\
&\quad \left. + \sum_{k=l+1}^n d_k \left\| \left[\mathbf{I}_{m \times m} + \boldsymbol{\Sigma}_{(-i)} \boldsymbol{\Sigma}_l^{-1} - \boldsymbol{\Sigma}_{(-l)} \boldsymbol{\Sigma}_{\mathbf{c}_l}^{-1} \right]^{-1} \right\|_2 \right\| \|\boldsymbol{\epsilon}_k^{j-1}\| \\
&< \delta \max_l \|\boldsymbol{\epsilon}_l^{j-1}\| \left\| \sum_{k=1, k \neq l}^n d_k \left\| \left[(\mathbf{I}_{m \times m} + \boldsymbol{\Sigma}_{(-l)} \boldsymbol{\Sigma}_i^{-1} - \boldsymbol{\Sigma}_{(-l)} \boldsymbol{\Sigma}_{\mathbf{c}_l}^{-1})^{-1} \right] \right\|_2 \right\| \\
&\leq \delta \max_l \|\boldsymbol{\epsilon}_l^{j-1}\|.
\end{aligned}$$

The result in (33) implies that $\max_i \|\boldsymbol{\epsilon}_i^j\| \leq \delta^j \max_i \|\boldsymbol{\epsilon}_i^0\|$ and $\max_i \|\boldsymbol{\epsilon}_i^j\| \rightarrow 0$ as $j \rightarrow \infty$. Therefore, \mathbf{z}_i^j 's converge to the values that satisfy equations (23) simultaneously. According to the result in Yee et al. (2002), \mathbf{z}_i^j converge to the posterior mean $E(\mathbf{z}_i | \mathbf{y}, \hat{\boldsymbol{\theta}}^{(k)})$ when $m \rightarrow \infty$. \diamond

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